Nice Banach Modules and Invariant Subspaces

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Abstract

Let \mathbb{A} be a semisimple unital commutative Banach algebra. We say that a Banach \mathbb{A} -module M is *nice* if every proper closed submodule of M is contained in a closed submodule of M of codimension 1. We provide examples of nice and non-nice modules.

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1 Introduction

In this article, all vector spaces are assumed to be over the field \mathbb{C} of complex numbers. As usual, \mathbb{R} is the field of real numbers, \mathbb{N} is the set of all positive integers, \mathbb{Z} is the set of integers and \mathbb{Z}_+ is the set of non-negative integers. For a Banach space X, L(X) stands for the algebra of bounded linear operators on X, while X^* is the space of continuous linear functionals on X. For $T \in L(X)$, its dual is denoted T^* : $T^* \in L(X^*)$, T'f(x) = f(Tx) for every $f \in X^*$ and every $x \in X$.

Throughout this article \mathbb{A} stands for a unital commutative semisimple Banach algebra. It is well-known and is a straightforward application of the Gelfand theory [2, 1] that for an ideal J in \mathbb{A} ,

$$J=\mathbb{A}\iff J\text{ is dense in }\mathbb{A}\iff\varkappa\big|_J\neq0\text{ for every }\varkappa\in\Omega(\mathbb{A}),$$

where $\Omega(\mathbb{A})$ is the spectrum of \mathbb{A} , that is, $\Omega(\mathbb{A})$ is the set of all (automatically continuous) non-zero algebra homomorphisms from \mathbb{A} to \mathbb{C} (endowed with the *-weak topology). Equivalently, every proper ideal in \mathbb{A} is contained in a closed ideal of codimension 1.

Let $\Omega^+(\mathbb{A})$ be the set of all algebra homomorphisms from \mathbb{A} to \mathbb{C} . That is, $\Omega^+(\mathbb{A})$ is $\Omega(\mathbb{A})$ together with the identically zero map from \mathbb{A} to \mathbb{C} . The main purpose of this paper is to draw attention to possible extensions of the above fact to Banach \mathbb{A} -modules. Clearly, each $\mathbb{A} \in \Omega^+(\mathbb{A})$ gives rise to the 1-dimensional \mathbb{A} -module $\mathbb{C}_{\mathbb{A}}$ being \mathbb{C} with the \mathbb{A} -module structure given by the multiplication $a\lambda = \mathbb{A}(a)\lambda$ for every $a \in \mathbb{A}$ and $\lambda \in \mathbb{C}$. It is also rather obvious that we have just listed all the 1-dimensional \mathbb{A} -modules up to an isomorphism.

Definition 1.1. Let M be a Banach \mathbb{A} -module. A *character* on M is a non-zero $\varphi \in M^*$ such that there exists $\varkappa \in \Omega^+(\mathbb{A})$ making φ into an \mathbb{A} -module morphism from M to \mathbb{C}_\varkappa .

Obviously, the kernel of a character on a Banach \mathbb{A} -module M is a closed \mathbb{A} -submodule of M.

Definition 1.2. Let M be a Banach \mathbb{A} -module. We say that M is *nice* if for every proper closed submodule of M is contained in a closed submodule of codimension 1. Equivalently, M is nice if and only if for every proper closed submodule N of M, there is a character φ on M such that φ vanishes on N.

The general question we would like to raise is:

Question 1.3. Characterize nice Banach A-modules.

The remark we started with ensures that \mathbb{A} is nice as an \mathbb{A} -module. In this paper we just present examples of nice and non-nice modules. Before even formulating the results, I would like to put forth my personal motivation for even looking at this question. Assume for a minute that \mathbb{A} is a subalgebra of L(X) for some Banach space X. We allow the norm topology of \mathbb{A} to be stronger (not necessarily strictly) than

the topology defined by the norm inherited from L(X). The multiplication $(A, x) \mapsto Ax$ defines a Banach \mathbb{A} -module structure on X. What are the characters on X? Why, one easily sees that they are exactly the common eigenvectors of A^* for $A \in \mathbb{A}$. What are the \mathbb{A} -submodules of X? They are exactly the invariant subspaces for the action of \mathbb{A} on X. Thus the \mathbb{A} -module X is nice exactly when every non-trivial closed \mathbb{A} -invariant subspace of X is contained in a closed \mathbb{A} -invariant hyperplane. Thus X being a nice \mathbb{A} -module translates into a strong and important property of the lattice of \mathbb{A} -invariant subspaces. Note that under relatively mild extra assumptions on \mathbb{A} , the nicety of X results in every closed \mathbb{A} -invariant subspace being the intersection of a collection of characters on X thus providing a complete description of the lattice of \mathbb{A} . A byproduct of this observation is the following easy example of a non-nice module.

Example 1.4. Let Ω be a non-empty compact subset of $\mathbb C$ with no isolated points and μ be a finite σ -additive purely non-atomic Borel measure on $\mathbb C$, whose support is exactly Ω . The pointwise multiplication equips $L^2(\Omega,\mu)$ with the structure of a Banach $C(\Omega)$ -module. This module is non-nice.

Proof. The $C(\Omega)$ -module $L^2(\Omega, \mu)$ does have plenty of closed submodules. For instance, every Borel subset A of Ω satisfying $\mu(A) \neq 0$ and $\mu(\Omega \setminus A) \neq 0$ generates a closed non-trivial submodule $M_A = \{f \in L^2(\Omega, \mu) : f \text{ vanishes outside } A\}$. On the other hand, we can always pick $f \in C(\Omega)$ satisfying $\mu(f^{-1}(\lambda)) = 0$ for every $\lambda \in \mathbb{C}$. In this case the dual of the multiplication by f operator on $L^2(\Omega, \mu)$ has empty point spectrum. Due to the above remark, our module possesses no characters at all (while possessing non-trivial closed submodules) and therefore can not possibly be nice.

In the positive direction we have the following two rather easy statements.

Proposition 1.5. The finitely generated free \mathbb{A} -module \mathbb{A}^n is nice.

Proposition 1.6. Let Ω be a Hausdorff compact topological space and X be a Banach space. Then the $C(\Omega)$ -module $C(\Omega, X)$ is nice, where $C(\Omega, X)$ carries the natural norm $||f|| = \sup\{||f(\omega)||_X : \omega \in \Omega\}$ and the module structure is given by the pointwise multiplication.

Note that Example 1.4 is rather cheatish since the non-nicety comes from the lack of characters. A really interesting situation is when a non-nice module possesses a separating set of characters. The following result says that this is quite possible. Recall that the Sobolev space $W^{1,2}[0,1]$ consists of the functions $f:[0,1]\to\mathbb{C}$ absolutely continuous on any bounded subinterval of I and such that $f'\in L_2[0,1]$. The space $W^{1,2}[0,1]$ with the inner product

$$\langle f, g \rangle_{1,2} = \int_0^1 (f(t)\overline{g(t)} + f'(t)\overline{g'(t)}) dt$$

is a separable Hilbert space. We denote $||f||_{1,2} = \sqrt{\langle f, f \rangle_{1,2}}$. Apart from being a Hilbert space, $W^{1,2}[0,1]$ is also a Banach algebra with respect to the pointwise multiplication (if one strives for the submultiplicativity of the norm together with the identity ||1|| = 1, he or she has to pass to an equivalent norm).

We say that a function f defined on [0,1] and taking values in a Banach space X is absolutely continuous if there exists an (automatically unique up to a Lebesgue-null set) Borel measurable function $g:[0,1]\to X$ such that

$$\int_{0}^{1} \|g(t)\| dt < +\infty \text{ and } \int_{0}^{x} g(t) dt = f(x) \text{ for each } x \in [0, 1],$$

where the second integral is considered in the Bochner sense. We denote the function g as f'. If H is a Hilbert space. The symbol $W^{1,2}([0,1],H)$ stands for the space of absolutely continuous functions $f:[0,1] \to H$ such that

$$\int_0^1 \|f'(t)\|^2 \, dt < +\infty.$$

The space $W^{1,2}([0,1],H)$ with the inner product

$$\langle f, g \rangle = \int_0^1 (\langle f(t), g(t) \rangle_H + \langle f'(t), g'(t) \rangle_H) dt$$

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is a Hilbert space and is separable if H is separable. In any case if $\{e_{\alpha}\}_{{\alpha}\in A}$ is an orthonormal basis of H, then the space $W^{1,2}([0,1],H)$ is naturally identified with the Hilbert direct sum of |A| copies of $W^{1,2}[0,1]$: $f\mapsto \{f_{\alpha}\}_{{\alpha}\in A}$, where $f_{\alpha}(t)=\langle f(t),e_{\alpha}\rangle_{H}$. It is also clear that $W^{1,2}([0,1],H)$ is naturally isomorphic to the Hilbert space tensor product of $W^{1,2}[0,1]$ and H. Clearly, $W^{1,2}([0,1],H)$ is a Banach $W^{1,2}[0,1]$ -module. This module possesses a lot of characters. Indeed, if $t\in [0,1]$ and $x\in H$, then the functional $f\mapsto \langle f(t),x\rangle_{H}$ is a character on $W^{1,2}([0,1],H)$. Moreover, these characters do separate points of $W^{1,2}([0,1],H)$.

Theorem 1.7. Let H be a Hilbert space. Then the $W^{1,2}[0,1]$ -module $W^{1,2}([0,1],H)$ is nice if and only if H is finite dimensional.

2 Proof of Proposition 1.6

It is easy to see that a character on $C(\Omega, X)$ is exactly a functional of the form

$$\varkappa_{\omega,\varphi}(f) = \varphi(f(\omega)), \text{ where } \omega \in \Omega \text{ and } \varphi \in X^* \setminus \{0\}.$$
(2.1)

The following lemma describes all closed submodules of $C(\Omega, X)$.

Lemma 2.1. Let M be a $C(\Omega)$ -submodule of $C(\Omega, X)$ and for each $\omega \in \Omega$ let $M_{\omega} = \{f(\omega) : f \in M\}$. Then the closure \overline{M} of M in $C(\Omega, X)$ satisfies

$$\overline{M} = \widetilde{M}, \quad where \ \widetilde{M} = \{ f \in C(\Omega, X) : f(\omega) \in \overline{M}_{\omega} \ for \ each \ \omega \in \Omega \},$$
 (2.2)

with \overline{M}_{ω} being the closure in X of M_{ω} .

Proof. Since $M \subseteq \widetilde{M}$ and \widetilde{M} is closed, we have $\overline{M} \subseteq \widetilde{M}$. Let $f \in \widetilde{M}$ and $\varepsilon > 0$. The desired equality will be verified if we show that there is $g \in M$ such that $||f - g|| < \varepsilon$. Indeed, in this case $\widetilde{M} \subseteq \overline{M}$ and therefore $\overline{M} = \widetilde{M}$.

Take $\omega \in \Omega$. Since M_{ω} is dense in \overline{M}_{ω} , there is $g_{\omega} \in M$ such that $||f(\omega) - g_{\omega}(\omega)||_{X} < \varepsilon$. Then $V_{\omega} = \{s \in \Omega : ||f(s) - g_{\omega}(s)||_{X} < \varepsilon\}$ is an open subset of Ω containing ω . Thus $\{V_{\omega}\}_{\omega \in \Omega}$ is an open covering of Ω . Since for every open covering of a Hausdorff compact topological space, there is a finite partition of unity consisting of continuous functions and subordinate to the covering [4], there are $\omega_1, \ldots, \omega_n \in \Omega$ and $\rho_1, \ldots, \rho_n \in C(\Omega)$ such that

$$0 \le \rho_j(s) \le 1$$
 for every $1 \le j \le n$ and $s \in \Omega$;
 $\rho_j(s) = 0$ whenever $1 \le j \le n$ and $s \in \Omega \setminus V_{\omega_j}$;
 $\rho_1(s) + \ldots + \rho_n(s) = 1$ for each $s \in \Omega$. (2.3)

Now we set $g = \rho_1 g_{\omega_1} + \ldots + \rho_n g_{\omega_n}$. Since M is a $C(\Omega)$ -module and $g_{\omega} \in M$, we have $g \in M$. Using (2.3) together with the inequality $||f(s) - g_{\omega_j}(s)||_X < \varepsilon$ for $s \in V_{\omega_j}$, we easily see that $||f(s) - g(s)||_X < \varepsilon$ for each $s \in \Omega$. Hence $g \in M$ and $||f - g|| < \varepsilon$, which completes the proof.

We are ready to prove Proposition 1.6. Let M be a closed submodule of $C(\Omega, X)$ such that none of the characters on $C(\Omega, X)$ vanishes on M. According to (2.1), the latter means that every $M_{\omega} = \{f(\omega) : f \in M\}$ is dense in X and therefore $\overline{M}_{\omega} = X$ for each $\omega \in \Omega$. Since M is closed, Lemma 2.1 says that $M = C(\Omega, X)$. The proof is complete.

3 Proof of Propositions 1.5

We start with the following easy observation. Let $\varkappa \in \Omega(\mathbb{A})$. Then the \mathbb{A} -module morphisms $\psi : \mathbb{A}^n \to \mathbb{C}_{\varkappa}$ are all given by

$$\varphi_c(a_1,\ldots,a_n) = \sum_{j=1}^n c_j \varkappa(a_j), \text{ where } c \in \mathbb{C}^n.$$

We shall prove a statement slightly stronger than Proposition 1.5.

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Proposition 3.1. Let $n \in \mathbb{N}$ and M be an \mathbb{A} -submodule of the free \mathbb{A} -module \mathbb{A}^n . Assume also that none of the characters on \mathbb{A}^n vanishes on M. Then $M = \mathbb{A}^n$.

Proof. We use induction with respect to n. The case n=1 is trivial (see the remark at the very start of the article). Assume now that $n \ge 2$ and that the conclusion of Proposition 1.5 holds for every smaller n. We interpret \mathbb{A}^n as $\mathbb{A}^n = \mathbb{A} \times \mathbb{A}^{n-1}$. The induction hypothesis easily implies that that the projection of M onto \mathbb{A}^{n-1} is onto. Let $J \subseteq \mathbb{A}$ be defined by $M \cap (\mathbb{A} \times \{0\}) = J \times \{0\}$. Then J is an ideal in A. If $J = \mathbb{A}$, we can factor out the first component in the product $\mathbb{A} \times \mathbb{A}^{n-1} = \mathbb{A}^n$ and then use the induction hypothesis to conclude that $M = \mathbb{A}^n$. Thus it remains to consider the case $J \neq \mathbb{A}$. Then there is $\mathbb{X} \in \Omega(\mathbb{A})$ such that $J \subseteq \ker \mathbb{X}$. Using the definition of J, and the facts that M is an \mathbb{A} -module, M projects onto the entire \mathbb{A}^{n-1} and \mathbb{X} vanishes on J, we can define $\psi : \mathbb{A}^{n-1} \to \mathbb{C}$ by the rule $\psi(b) = \mathbb{X}(a)$ if $(a,b) \in M \subseteq \mathbb{A} \times \mathbb{A}^{n-1}$. It is easy to see that ψ is a well-defined continuous linear functional and that $\psi : \mathbb{A}^{n-1} \to \mathbb{C}_{\mathbb{X}}$ is an \mathbb{A} -module morphism. According to the above display there are $c_1, \ldots, c_{n-1} \in \mathbb{C}$ such that $\psi(a_1, \ldots, a_{n-1}) = \sum_{j=1}^{n-1} c_j \mathbb{X}(a_j)$ for every $a_1, \ldots, a_{n-1} \in \mathbb{A}$. By definition of ψ , we now see that

 $\varphi: \mathbb{A}^n \to \mathbb{C}$ vanishes on M, where φ is defined by the formula $\varphi(a_1, \ldots, a_n) = \sum_{j=1}^n c_j \varkappa(a_j)$ with $c_n = -1$.

By the above display, $\varphi : \mathbb{A}^n \to \mathbb{C}_{\varkappa}$ is an \mathbb{A} -module morphism. Since $c_n \neq 0$, $\varphi \neq 0$ and therefore φ is a character on \mathbb{A}^n . We have produced a character on \mathbb{A}^n vanishing on M, which contradicts the assumptions. Thus the case $J \neq \mathbb{A}$ does not occur, which completes the proof.

4 Proof of Theorem 1.7

In this section, for a function f on an interval I of the real line $||f||_2$ will always denote the L^2 -norm of f (with respect to the Lebesgue measure), while $||f||_{\infty}$ always stands for the L^{∞} -norm of f.

Lemma 4.1. Let $-\infty < \alpha < \beta < +\infty$, $a,b \in \mathbb{C}$ and $\varepsilon > 0$. Then there exists $f \in C^1[\alpha,\beta]$ such that $f(\alpha) = f(\beta) = 0$, $f'(\alpha) = a$, $f'(\beta) = b$ and $||f||_{\infty} < \varepsilon$.

Proof. Let $\varphi \in C^1[0,\infty)$ be a monotonically non-increasing function such that $\varphi(0)=1,\ \varphi'(0)=0$ and $\varphi(x)=0$ for $x\geqslant 1$. For any $\delta\in(0,\frac{\beta-\alpha}{2})$ let

$$f_{\delta}(x) = \begin{cases} 0 & \text{if } x \in (\alpha + \delta, \beta - \delta), \\ a(x - \alpha)\varphi((x - \alpha)/\delta) & \text{if } x \in [\alpha, \alpha + \delta), \\ b(x - \beta)\varphi((\beta - x)/\delta) & \text{if } x \in (\beta - \delta, \beta]. \end{cases}$$

Obviously, $f_{\delta} \in C^1[\alpha, \beta]$, $f_{\delta}(\alpha) = f_{\delta}(\beta) = 0$, $f'_{\delta}(\alpha) = a$, $f'_{\delta}(\beta) = b$ and $||f||_{\infty} \leq \delta \max\{|a|, |b|\}$. Hence the function $f = f_{\delta}$ for $\delta < \varepsilon / \max\{|a|, |b|\}$ satisfies all desired conditions.

Lemma 4.2. Let $K \subset [0,1]$ be a nowhere dense compact set, $a \in C(K)$, $f \in C[0,1]$ and $\varepsilon > 0$. Then there exists $g \in C^1[0,1]$ such that $g'|_K = a$ and $\|g - f\|_{\infty} < \varepsilon$.

Proof. Since $C^1[0,1]$ is dense in the Banach space C[0,1], we can, without loss of generality, assume that $f \in C^1[0,1]$. Since any continuous function on K admits a continuous extension to [0,1] (one can apply, for instance, the Tietze theorem [4]), there exists $h \in C[0,1]$ such that h(x) = a(x) - f'(x) for any $x \in K$. Let $\delta > 0$. Since K is nowhere dense, there exist

$$0 = \alpha_1 < \beta_1 < \alpha_2 < \beta_2 < \dots < \alpha_n < \beta_n = 1$$

such that $\beta_j - \alpha_j < \varepsilon$ for any j = 1, ..., n and $K \subset \bigcup_{j=1}^n I_j$, where $I_j = [\alpha_j, \beta_j]$. Let

$$a_j = \int_{\alpha_j}^{\beta_j} h(t) dt$$
 for $1 \le j \le n - 1$.

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By Lemma 4.1, for $1 \leqslant j \leqslant n-1$, there is $\varphi_j \in C^1[\beta_j, \alpha_{j+1}]$ such that $\varphi_j(\beta_j) = \varphi_j(\alpha_{j+1}) = 0$, $\varphi_j'(\beta_j) = h(\beta_j) + \frac{a_j}{\alpha_{j+1} - \beta_j}$, $\varphi_j'(\alpha_{j+1}) = h(\alpha_{j+1}) + \frac{a_j}{\alpha_{j+1} - \beta_j}$ and $\|\varphi_j\|_{\infty} < \delta$. Consider the function

$$\psi(x) = \begin{cases} \int_{\alpha_j}^x h(t) dt & \text{if } x \in [\alpha_j, \beta_j], \ 1 \leqslant j \leqslant n, \\ \varphi_j(x) + \frac{a_j(x - \alpha_{j+1})}{\beta_j - \alpha_{j+1}} & \text{if } x \in (\beta_j, \alpha_{j+1}), \ 1 \leqslant j \leqslant n - 1. \end{cases}$$

The values of φ_j' at β_j and α_{j+1} were chosen in such a way that $\psi \in C^1[0,1]$. Moreover, $\psi'\big|_{I_j} = h$ for $1 \leqslant j \leqslant n$. Hence, $(\psi+f)'\big|_K = a$. Let us estimate $\|\psi\|_{\infty}$. If $1 \leqslant j \leqslant n-1$ and $x \in [\beta_j,\alpha_{j+1}]$, then $|\psi(x)| \leqslant \delta + |a_j| \leqslant \delta + |\beta_j - \alpha_j| \|h\|_{\infty} \leqslant \delta(1 + \|h\|_{\infty})$. If $1 \leqslant j \leqslant n$ and $x \in [\alpha_j,\beta_j]$, then $|\psi(x)| \leqslant |\beta_j - \alpha_j| \|h\|_{\infty} \leqslant \delta \|h\|_{\infty}$. Hence $\|\psi\|_{\infty} \leqslant \delta(1 + \|h\|_{\infty})$. Choose $\delta < \varepsilon/(1 + \|h\|_{\infty})$ and denote $g = \psi + f$. Then $g'\big|_K = a$ and $\|g - f\|_{\infty} = \|\psi\|_{\infty} < \varepsilon$.

Lemma 4.3. Let $K \subset [0,1]$ be a nowhere dense compact set and $\varepsilon > 0$. Then there exists $f \in C(K)$ such that

$$\int_{K} f(t) dt = 0 \quad and \quad \|\chi + g\|_{2} \leqslant \varepsilon, \quad where \quad g(x) = \int_{K \cap [x,1]} f(t) dt$$

and χ is the indicator function of K ($\chi(x) = 1$ if $x \in K$ and $\chi(x) = 0$ if $x \in [0,1] \setminus K$).

Proof. If the Lebesgue measure $\mu(K)$ of K is zero, the statement is trivially true since the function $f \equiv 0$ satisfies the desired conditions for any $\varepsilon > 0$. Thus, we can assume that $\mu(K) > 0$. Let $n \in \mathbb{N}$. Since K is nowhere dense and has positive Lebesgue measure, we can choose $n \in \mathbb{N}$ and $\alpha_k, \beta_k, a_k, b_k, u_k, v_k \in [0, 1] \setminus K$ for $1 \leq k \leq n$ in such a way that

$$\alpha_k < \beta_k < a_k < b_k < u_k < v_k \text{ for } 1 \leqslant k \leqslant n \text{ and } v_{k-1} < \alpha_k \text{ for } 2 \leqslant k \leqslant n,$$

$$0 < \mu(K \cap [\alpha_k, \beta_k]) < \frac{\varepsilon^2}{16n} \text{ and } 0 < \mu(K \cap [u_k, v_k]) < \frac{\varepsilon^2}{16n} \text{ for } 1 \leqslant k \leqslant n,$$

$$(4.1)$$

$$\mu\left(\left(\bigcup_{k=1}^{n} \left[\alpha_k, v_k\right]\right) \setminus K\right) < \frac{\varepsilon^2}{8}.\tag{4.2}$$

Consider the function $f: K \to \mathbb{R}$ defined by the formula

$$f(x) = \begin{cases} \frac{1}{\mu(K \cap [\alpha_k, \beta_k])} & \text{if } x \in K \cap [\alpha_k, \beta_k], \ 1 \leqslant k \leqslant n; \\ \frac{-1}{\mu(K \cap [u_k, v_k])} & \text{if } x \in K \cap [u_k, v_k], \ 1 \leqslant k \leqslant n; \\ 0 & \text{otherwise.} \end{cases}$$

Obviously $f \in C(K)$ and

$$\int_{K} f(t) dt = \sum_{k=1}^{n} \left(\int_{K \cap [\alpha_{k}, \beta_{k}]} f(t) dt - \int_{K \cap [u_{k}, v_{k}]} f(t) dt \right) = \sum_{k=1}^{n} (1 - 1) = 0.$$

Let $g:[0,1]\to\mathbb{R}$ be defined by

$$g(x) = \int_{K \cap [x,1]} f(t) dt.$$

From the definition of f it follows that $|g(x)| \leq 1$ for any $x \in [0,1]$, g(x) = -1 if $x \in \bigcup_{k=1}^{n} [\beta_k, u_k]$ and

 $\chi(x) = g(x) = 0$ if $x \in [0,1] \setminus \bigcup_{k=1}^{n} [\alpha_k, v_k]$. Hence the set $\Omega = \{x \in [0,1] : g(x) + \chi(x) \neq 0\}$ is contained in the union

$$\Omega_1 = \left(\left(\bigcup_{k=1}^n [\alpha_k, v_k] \right) \setminus K \right) \cup \left(\bigcup_{k=1}^n ([\alpha_k, \beta_k] \cap K) \right) \cup \left(\bigcup_{k=1}^n ([u_k, v_k] \cap K) \right).$$

Therefore

$$||g + \chi||_2^2 = \int_0^1 (g(x) + \chi(x))^2 \le 4\mu(\Omega) \le 4\mu(\Omega_1).$$

Using (4.1) and (4.2), we see that $\mu(\Omega_1) \leq \varepsilon^2/4$. Hence $||g + \chi||_2 \leq \varepsilon$.

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Lemma 4.4. Let $\{e_n\}_{n\in\mathbb{Z}_+}$ be an orthonormal basis in a separable Hilbert space H and scalar sequences $\{\gamma_n\}_{n\in\mathbb{N}}$ and $\{\delta_n\}_{n\in\mathbb{N}}$ be such that

$$\sum_{n=1}^{\infty} (|\gamma_n|^2 + |\delta_n|^2) < \infty. \tag{4.3}$$

Let also $f_0 = e_0 + \sum_{n=1}^{\infty} \gamma_n e_n$ and $f_n = e_n - \delta_n e_0$ for $n \in \mathbb{N}$. Then the linear span of $\{f_n : n \in \mathbb{Z}_+\}$ is dense in H if and only if

$$\sum_{n=1}^{\infty} \gamma_n \delta_n \neq -1. \tag{4.4}$$

Proof. Condition (4.3) implies that the linear operator $T: H \to H$ such that $Te_0 = \sum_{n=1}^{\infty} \gamma_n e_n$ and $Te_n = -\delta_n e_0$ for $n \in \mathbb{N}$ is bounded. Since the range of T is at most two-dimensional, T is compact. By the Fredholm theorem [3], the operator S = I + T has dense range if and only if S is injective. Since $Se_n = f_n$ for $n \in \mathbb{Z}_+$, the linear span of $\{f_n\}_{n \in \mathbb{Z}_+}$ is dense in H if and only if the operator S injective.

The equation Sx = 0, $x \in H$ can be rewritten as

$$\langle x, e_0 \rangle \left(1 + \sum_{n=1}^{\infty} \gamma_n \delta_n \right) = 0 \text{ and } \langle x, e_n \rangle = \gamma_n \langle x, e_0 \rangle \text{ for any } n \in \mathbb{N}.$$

If $\sum_{n=1}^{\infty} \gamma_n \delta_n \neq -1$, the first equation implies $\langle x, e_0 \rangle = 0$ and the rest yield $\langle x, e_n \rangle = 0$ for each $n \in \mathbb{N}$. Thus in this case x = 0. That is, S is injective and therefore the linear span of $\{f_n : n \in \mathbb{Z}_+\}$ is dense in H. If $\sum_{n=1}^{\infty} \gamma_n \delta_n = -1$, the system of the equations in the above display has the non-zero solution $x = x_0 + \sum_{n=1}^{\infty} \gamma_n e_n \in H$. Hence S is not injective and therefore the linear span of $\{f_n : n \in \mathbb{Z}_+\}$ is non-dense.

We are ready to prove Theorem 1.7. First, if $n \in \mathbb{N}$ and H is n-dimensional, then $W^{1,2}([0,1], H)$ is isomorphic to the free $W^{1,2}[0,1]$ -module with n generators and the nicety of $W^{1,2}([0,1], H)$ follows from Proposition 1.5. It is easy to see that a direct (module) summand of a nice module is nice. Thus the proof of Theorem 1.7 will be complete if we verify that $W^{1,2}([0,1],\ell_2)$ is non-nice. In order to do this, we have to construct a proper closed $W^{1,2}[0,1]$ -submodule M of $W^{1,2}([0,1],\ell_2)$ such that none of the characters on $W^{1,2}([0,1],\ell_2)$ vanishes on M. Now we shall do just that.

Pick a nowhere dense compact set $K \subset [0,1]$ of positive Lebesgue measure and let χ be the indicator function of K. By Lemma 4.3, there exists $A_n \in C(K)$ such that for any $n \in \mathbb{N}$,

$$\int_{K} A_n(x) dx = 0, \tag{4.5}$$

$$||B_n + \chi||_2 < 2^{-n}$$
, where $B_n(x) = \int_{K \cap [x,1]} A(t) dt$. (4.6)

We also set $A_0 = 0$, $B_0 = 0$ and $S_0 = 1$. By Lemma 4.2, there exist $S_n \in C^1[0,1]$ such that

$$S'_n|_K = A_n \text{ and } ||S_n - 1||_{\infty} < 2^{-n} \text{ for each } n \in \mathbb{N}.$$
 (4.7)

Denote $\rho_n = n^2(S_n - S_{n-1})$ for $n \in \mathbb{N}$. Then $\rho_n \in C^1[0,1]$ and according to (4.7),

$$\|\rho_n\|_{\infty} \le n^2 (\|S_n - 1\|_{\infty} + \|S_{n-1} - 1\|_{\infty}) \le n^2 (2^{1-n} + 2^{-n}) = 3n^2 2^{-n} \text{ for each } n \in \mathbb{N}.$$
 (4.8)

Let also $\{e_n\}_{n\in\mathbb{Z}_+}$ be the standard orthonormal basis in ℓ_2 . Consider the functions $f^{[n]}\in W^1_2([0,1],\ell_2)$ defined by the formulas

$$f^{[0]}(x) = e_0 + \sum_{n=1}^{\infty} n^{-2} e_n$$
 and $f^{[n]}(x) = e_n - \rho_n(x) e_0$ for $n \in \mathbb{N}$.

Let now M be the closed $W^{1,2}[0,1]$ -submodule of $W^{1,2}([0,1],\ell_2)$ generated by the set $\{f^{[n]}:n\in\mathbb{Z}_+\}$. Equivalently, M is the closed linear span in $W^{1,2}([0,1],\ell_2)$ of the set $\{\varphi f^{[n]}:n\in\mathbb{Z}_+,\ \varphi\in W^{1,2}[0,1]\}$. It is easy to see that every character on $W^{1,2}([0,1],\ell_2)$ has the shape

$$\varphi_{t,y}(f) = \langle f(t), y \rangle_H$$
, where $t \in [0, 1]$ and $y \in \ell_2 \setminus \{0\}$.

Thus in order for every character on $W^{1,2}([0,1], \ell_2)$ not to vanish on M it is necessary and sufficient for $M_t = \{f(t) : f \in M\}$ to be dense in ℓ_2 for every $t \in [0,1]$. Let $t \in [0,1]$. By definition of ρ_n and (4.7), we have

$$\sum_{n=1}^{\infty} n^{-2} \rho_n(t) = \lim_{m \to \infty} \sum_{n=1}^{m} (S_n(t) - S_{n-1}(t)) = \lim_{m \to \infty} (S_m(t) - S_0(t)) = 0 \neq -1.$$
 (4.9)

By Lemma 4.4 with $\gamma_n = n^{-2}$ and $\delta_n = \rho_n(t)$, the linear span of $\{f^{[n]}(t)\}_{n \in \mathbb{Z}_+}$ is dense in ℓ_2 . Since $f^{[n]} \in M$, M_t is dense in ℓ_2 . Thus none of the characters on $W^{1,2}([0,1],\ell_2)$ vanishes on M. It remains to verify that $M \neq W^{1,2}([0,1],\ell_2)$. Consider $g_n \in W_2^1[0,1]^*$ for $n \in \mathbb{Z}_+$, defined by the formula

$$g_n(\varphi) = \int_K (\rho_n \varphi)'(x) dx$$
, where ρ_0 is assumed to be identically 1.

We start with estimating the norms of the functionals g_n . Clearly,

$$g_n(\varphi) = \int_K \rho_n(x)\varphi'(x) dx + \int_K \rho'_n(x)\varphi(x) dx \text{ for any } \varphi \in W^{1,2}[0,1].$$
 (4.10)

Since $\rho'_n(x) = n^2(S'_n(x) - S'_{n-1}(x)) = n^2(A_n(x) - A_{n-1}(x))$ for $x \in K$, we have

$$\int_{K} \rho'_{n}(x)\varphi(x) dx = n^{2} \int_{K} (A_{n}(x) - A_{n-1}(x))\varphi(x) dx = n^{2} \int_{0}^{1} (B'_{n-1}(x) - B'_{n}(x))\varphi(x) dx.$$

By (4.5) and (4.6),

$$B_n(0) = B_n(1) = 0$$
 for $n \in \mathbb{Z}_+$.

Integrating by parts and using the above display, we obtain

$$\int_K \rho'_n(x)\varphi(x) \, dx = n^2 \int_0^1 (B'_{n-1}(x) - B'_n(x))\varphi(x) \, dx = n^2 \int_0^1 (B_n(x) - B_{n-1}(x))\varphi'(x) \, dx.$$

This formula together with (4.10) yields

$$|g_n(\varphi)| \le ||\varphi'||_2 (||\rho_n||_2 + n^2 ||B_n - B_{n-1}||_2)$$
 for $n \in \mathbb{N}$.

Since $\|\rho_n\|_2 \leqslant 3n^2 2^{-n}$ and $\|B_n - B_{n-1}\|_2 \leqslant \|B_n + \chi\|_2 + \|B_{n-1} + \chi\|_2 \leqslant 2^{1-n} + 2^{-n} = 3 \cdot 2^{-n}$, we have $|g_n(\varphi)| \leqslant 6n^2 2^{-n} \|\varphi\|_{1,2}$. Hence $\|g_n\| \leqslant 6n^2 2^{-n}$ for each $n \in \mathbb{N}$. Therefore $\sum_{n=0}^{\infty} \|g_n\|^2 < \infty$. Thus the formula

$$g(h) = \sum_{n=0}^{\infty} g_n(h_n)$$

defines a continuous linear functional on $W^{1,2}([0,1],\ell_2)$, where, as usual, $h_n(t) = \langle h(t), e_n \rangle$. Since $g_0 \neq 0$, we have $g \neq 0$. In order to show that $M \neq W^{1,2}([0,1],\ell_2)$, it suffices to verify that g(h) = 0 for any $h \in M$.

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For this it is enough to check that $g(\varphi f^{[n]}) = 0$ for every $\varphi \in W^{1,2}[0,1]$ and $n \in \mathbb{Z}_+$. First, let $n \in \mathbb{N}$. Then by definition of g_n , we immediately have

$$g(\varphi f^{[n]}) = g_n(\varphi) - g_0(\rho_n \varphi) = 0.$$

It remains to prove that $g(\varphi f^{[0]}) = 0$. Using the uniform convergence of the series $\sum_{n=1}^{\infty} n^{-2} \rho_n$ provided by the estimate (4.8), we have

$$g(\varphi f^{[0]}) = g_0(\varphi) + \sum_{n=1}^{\infty} n^{-2} g_n(\varphi) = \int_K \left(\varphi'(x) + \sum_{n=1}^{\infty} n^{-2} (\rho_n \varphi)'(x) \right) dx =$$

$$= \int_K \varphi'(x) \left(1 + \sum_{n=1}^{\infty} n^{-2} \rho_n(x) \right) dx + \lim_{m \to \infty} \int_K \varphi(x) \left(\sum_{n=1}^m n^{-2} \rho'_n(x) \right) dx.$$

By (4.9), $\sum_{n=1}^{\infty} n^{-2} \rho_n(x) \equiv 0$. On the other hand, using (4.7) and the equality $S_0 = 1$, we have

$$\sum_{n=1}^{m} n^{-2} \rho'_n(x) = \sum_{n=1}^{m} (S'_n(x) - S'_{n-1}(x)) = S'_m(x) = A_m(x) \text{ for each } x \in K.$$

Hence

$$g(\varphi f^{[0]}) = \int_{K} \varphi'(x) dx + \lim_{m \to \infty} \int_{K} \varphi(x) A_m(x) dx.$$

$$(4.11)$$

Integrating by parts, we obtain

$$\int_{K} \varphi(x) A_{m}(x) dx = -\int_{0}^{1} \varphi(x) B'_{m}(x) dx = \int_{0}^{1} \varphi'(x) B_{m}(x) dx = \int_{0}^{1} \varphi'(x) (B_{m}(x) + \chi(x)) dx - \int_{K} \varphi'(x) dx.$$

According to (4.11) and the above display,

$$g(\varphi f^{[0]}) = \lim_{m \to \infty} \int_{0}^{1} \varphi'(x) (B_m(x) + \chi(x)) dx.$$
 (4.12)

By (4.6) and (4.12), $g(\varphi f^{[0]}) = 0$ for each $\varphi \in W^{1,2}[0,1]$. Thus g(h) = 0 for every $h \in M$ and therefore $M \neq W^{1,2}([0,1], \ell_2)$. The proof of Theorem 1.7 is complete.

5 Remarks

One can easily generalize Theorem 1.7 by taking most any algebra of smooth functions instead of $W^{1,2}[0,1]$. For example, following the same route of argument with few appropriate amendments one can show that if X is an infinite dimensional Banach space and $k \in \mathbb{N}$, then $C^k([0,1],X)$ as a $C^k[0,1]$ -module is non-nice. We opted for $W^{1,2}([0,1],H)$ to make a point that even the friendly Hilbert space environment does not save the day.

Theorem 1.7 says that there are weird proper closed submodules of $W^{1,2}([0,1], \ell_2)$ which are not contained in any closed submodule of codimension 1. The following question remains wide open.

Question 5.1. Characterize closed submodules of $W^{1,2}([0,1],\ell_2)$.

References

- [1] H. Dales, Banach algebras and automatic continuity, Oxford University Press, New York, 2000
- [2] A. Helemsky, Banach and locally convex algebras, Oxford University Press, New York, 1993
- [3] W. Rudin Functional Analysis, McGra2-Hill, New York, 1991
- [4] R. Engelking, General Topology, Heldermann, Berlin, 1989